# A continuous Gaussian approximation to a nonparametric regression in two dimensions 

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Estimating the mean in a nonparametric regression on a two-dimensional regular grid of design points is asymptotically equivalent to estimating the drift of a continuous Gaussian process on the unit square. In particular, we provide a construction of a Brownian sheet process with a drift that is almost the mean function in the nonparametric regression. This can be used to apply estimation or testing procedures from the continuous process to the regression experiment as in Le Cam's theory of equivalent experiments. Our result is motivated by first looking at the amount of information lost in binning the data in a density estimation problem.

Keywords: asymptotic equivalence of experiments; density estimation; nonparametric regression

## 1. Introduction

The purpose of this paper is to establish a connection between a nonparametric regression on a two-dimensional set of design points and an appropriate continuous Gaussian approximation. This connection provides a bound on Le Cam's deficiency distance between the experiments and allows inference in the easier problem (the continuous case) to be applied to the practical problem (observations on a finite grid of points). The motivation for the form of this connection comes from a similar result: approximating the problem of estimating an unknown density from $n$ independent observations by the experiment that observes only those observations aggregated into $m$ bins.

Brown and Low (1996) showed that the nonparametric regression experiment

$$
\begin{equation*}
Y_{i}=g\left(x_{i}\right)+\sigma \xi_{i}, \quad i=1, \ldots, n \tag{1}
\end{equation*}
$$

with $\xi_{i}$ independent standard normals, $x_{i}=i /(n+1), \sigma$ a known value, and $g$ an unknown smooth function on $[0,1]$, is asymptotically equivalent to an observation of a continuous Gaussian process with an unknown drift function

$$
\begin{equation*}
\mathrm{d} Y_{t}=g(t) \mathrm{d} t+\frac{\sigma}{\sqrt{n}} \mathrm{~d} W_{t}, \tag{2}
\end{equation*}
$$

where $W_{t}$ is a standard Brownian motion on [0,1] whenever the class of possible mean functions $g \in \mathcal{G}$ is a subset of a Lipschitz $(\alpha)$ space with $\alpha>1 / 2$.

This result does not immediately extend to a regression on a two-dimensional space. The higher-dimensional result requires that the class of drift functions be smoother; in particular,
they must be differentiable, $\alpha>1$. We propose a construction that can take advantage of this added smoothness.

The idea of developing asymptotic results for nonparametric regression by appealing to a continuous Gaussian processes approximation is widely used; see, for example, Donoho et al. (1995) or Efromovich (1999, Chapter 7). Donoho and Johnstone (1999) described a method for constructing wavelet coefficients from the nonparametric regression as if the original process was continuous. They showed that the squared-error loss was not significantly affected by the approximation. Our construction leads to the same statistical estimators. Brown et al. (2002) extended their original result to include regression with a random design, but our result assumes a fixed grid of equally spaced design points. Grama and Nussbaum (1998) established equivalence for regression problems with non-normal errors.

Nussbaum (1996) showed that observing $n$ independent observations from an unknown density is equivalent to a Brownian motion plus drift. Carter (2002) and Brown et al. (2004) constructed a connection between these experiments by comparing their behaviour on a finite partition of the unit interval. It is also then necessary to bound the error in 'discretizing' the continuous observations. The process in (2) is approximated by its increments $Y(j / m)-Y((j-1) / m)$, and the independent observations are approximated by the number falling in each of $m$ equal subintervals. These two problems are related and the solution of the binning problem will motivate the approach used to go from the continuous to the finite-dimensional Gaussian process. The technical bounds in Section 4 can essentially be used in both situations.

### 1.1. Le Cam's deficiency distance

The constructions described in Sections 2 and 3 will bound the deficiency distance between the experiments. This statistical distance compares the relevant information about the parameter that is available in the two sets of distributions.

First, the total variation distance between two distributions is

$$
\begin{equation*}
\|P-Q\|=\sup _{A \in \mathcal{A}} P(A)-Q(A) \tag{3}
\end{equation*}
$$

where $P$ and $Q$ are both measures on the $\sigma$-field $\mathcal{A}$. This distance also bounds the difference in the expectations of bounded functions $|g|<1$,

$$
\sup _{\{g:|g| \leqslant 1\}}|P(g)-Q(g)| \leqslant 2\|P-Q\| .
$$

However, this distance is equal to 1 if $P$ and $Q$ do not have common support.
A statistical experiment $\mathcal{P}$ consists of a set of distributions $\left\{P_{\theta}: \theta \in \Theta\right\}$ indexed by a parameter set $\Theta$ for data $X \in \mathcal{X}$ with a $\sigma$-field $\mathcal{A}$. A second experiment $\mathcal{Q}=\left\{Q_{\theta}: \theta \in \Theta\right\}$ has the same parameter set but a different sample space $(\mathcal{Y}, \mathcal{B})$. To compare these two experiments, we need a way to connect the sample space $\mathcal{Y}$ to $\mathcal{X}$. A randomization of the data $X$ can be described by the conditional distribution $K_{x}$ on $(\mathcal{Y}, \mathcal{B})$ given $X=x$. Let $P_{\theta} K_{x}$ represent the marginal distribution on $(\mathcal{Y}, \mathcal{B})$. The randomization does not depend on
$\theta$ so $P_{\theta} K_{x}$ cannot contain any additional information about $\theta$. If $\left\|P_{\theta} K_{x}-Q_{\theta}\right\|$ is always small, then $\mathcal{Q}$ does not have much more information about $\theta$ than $\mathcal{P}$.

Le Cam's deficiency (Le Cam 1986, pp. 18-20) is

$$
\delta(\mathcal{P}, \mathcal{Q})=\inf _{K} \sup _{\theta \in \Theta}\left\|P_{\theta} K-Q_{\theta}\right\|
$$

where $K$ is a 'transition' from the linear space including probability distributions $P_{\theta}$ to a space that includes the measures $Q_{\theta}$. For experiments that depend on an increasing sample size $n$, if $\delta\left(\mathcal{P}_{n}, \mathcal{Q}_{n}\right) \rightarrow 0$ then the sequence of experiments $\mathcal{P}_{n}$ is asymptotically as informative as $\mathcal{Q}_{n}$. Furthermore, if $\max \left[\delta\left(\mathcal{P}_{n}, \mathcal{Q}_{n}\right), \delta\left(\mathcal{Q}_{n}, \mathcal{P}_{n}\right)\right] \rightarrow 0$ then $\mathcal{P}_{n}$ and $\mathcal{Q}_{n}$ are termed asymptotically equivalent.

For our purposes, it is not necessary to think in terms of these general transitions. We will bound the deficiency using a transformation of the observations $X$ from $\mathcal{P}$ that may also include an external randomization $T(X, W)$. The random variable $W$ represents the external randomization that has the same distribution for all $\theta$. Then $K_{x}$ is the conditional distribution of $T(X, W)$ given $X=x$. Thus,

$$
\begin{aligned}
\delta(\mathcal{P}, \mathcal{Q}) & \leqslant \sup _{\theta \in \Theta}\left\|P_{\theta} K_{x}-Q_{\theta}\right\| \\
& =\sup _{\theta \in \Theta} \sup _{A \in \mathcal{A}}\left|P_{\theta}\{T(X, W) \in A\}-Q_{\theta}\{Y \in A\}\right|
\end{aligned}
$$

The usefulness of a bound of this type is in its flexibility. Suppose that $\sup _{\theta \in \Theta}\left\|P_{\theta} K_{x}-Q_{\theta}\right\| \leqslant \tau$. Then for any bounded loss function $|L(a, \theta)| \leqslant 1$, any decision procedure $d(Y)$ in the experiment $\mathcal{Q}$ which has risk $R(d(Y), \theta)=Q_{\theta} L(d(Y), \theta)$ generates a randomized decision procedure $d[T(X, W)]$ such that

$$
R(d[T(X, W)], \theta)=P_{\theta} L(d[T(X, W)], \theta) \leqslant Q_{\theta} L(d(Y), \theta)+2 \tau
$$

In other words, the $T$ transformation maps good decision procedures for $\mathcal{Q}$ to good decision procedures for the $\mathcal{P}$ experiment.

### 1.2. Main results

The parameter sets in these nonparametric experiments are Lipschitz classes of differentiable functions $\mathcal{L}(\alpha, M)$ for $1<\alpha \leqslant 2$ where $f \in \mathcal{L}(\alpha, M)$ implies that $|f(x)| \leqslant M,\left|f^{\prime}(x)\right| \leqslant M$, and

$$
\begin{equation*}
\left|f^{\prime}(x)-f^{\prime}(y)\right| \leqslant M|x-y|^{\alpha-1} \tag{4}
\end{equation*}
$$

for every $x$ and $y$ in the sample space.
For $g:[0,1]^{2} \mapsto \mathbb{R}$, the analogous conditions hold with the Euclidean norm in $\mathbb{R}^{2}$ replacing absolute values. The partial derivatives exist and the vector of partials, $g^{\prime}(\mathbf{x})$, is bounded and smooth.

$$
\begin{aligned}
\mathcal{L}(\alpha, M)= & \left\{g: \sup _{\mathbf{x} \in[0,1]^{2}}|g(\mathbf{x})| \leqslant M, \sup _{\mathbf{x} \in[0,1]^{2}}\left|g^{\prime}(\mathbf{x})\right| \leqslant M,\right. \\
& \left.\sup _{\mathbf{x}, \mathbf{y} \in[0,1]^{2}}\left|g^{\prime}(\mathbf{x})-g^{\prime}(\mathbf{y})\right| \leqslant M|\mathbf{x}-\mathbf{y}|^{\alpha-1}\right\} .
\end{aligned}
$$

Theorem 1. For $n$ equally spaced design points $\mathbf{x}_{i, j}^{*}$ in $[0,1]^{2}$,

$$
\mathbf{x}_{i, j}^{*}=\left(\frac{2 i-1}{2 \sqrt{n}}, \frac{2 j-1}{2 \sqrt{n}}\right)
$$

let $\overline{\mathbb{Q}}_{g}$ be the distribution of the

$$
Y_{i, j}=g\left(\mathbf{x}_{i, j}^{*}\right)+\sigma \xi_{i, j}
$$

where $\sigma$ is known, the $\xi_{i, j}$ are independent standard normals, and $g$ is an unknown function from $[0,1]^{2}$ to $\mathbb{R}$ that is in $\mathcal{L}(\alpha, M)$. Let $\mathbb{Q}_{g}$ be the distribution of the Gaussian process

$$
\begin{equation*}
Y(\mathbf{t})=\int_{0}^{t_{1}} \int_{0}^{t_{2}} g(\mathbf{x}) \mathrm{d} x_{2} \mathrm{~d} x_{1}+\frac{\sigma}{\sqrt{n}} W(\mathbf{t}), \tag{5}
\end{equation*}
$$

where $W(\mathbf{t})$ is a Brownian sheet on the unit square. Then there exists a randomization $K_{y}$ such that

$$
\sup _{g \in \mathcal{L}(\alpha, M)}\left\|\overline{\mathbb{Q}}_{g} K_{y}-\mathbb{Q}_{g}\right\| \leqslant \frac{6 M}{\sqrt{2 \pi \sigma^{2}}}\left(n^{1 / 2-\alpha / 2}+n^{-1 / 4}\right)
$$

For $\alpha>1$, this implies that the error made by performing the inference in the continuous experiment is asymptotically negligible. This is a reasonably sharp result in that Brown and Zhang (1998) showed that the experiments are not equivalent for $\alpha=1$. Theorem 1 is proved in Section 3.

The experiments are asymptotically equivalent because there is a simple transformation in the other direction, taking the increments of the Gaussian process, that produces a smaller error than in Theorem 1. Details are given in Section 4.3.

It will be instructive to first show the following result about density estimation experiments. Let $\mathbb{P}_{f}=P_{f}^{n}$ be the distribution of $n$ independent observations from $P_{f}$, a distribution on the unit interval with density $f$. Then let $\mathbb{P}_{f}$ be the distribution that results from binning the $n$ observations into $m$ equal-length subintervals. The new observations have an $m$-dimensional multinomial distribution.

Theorem 2. For a density $f \in \mathcal{F}(\alpha, M, \varepsilon)$ such that $\mathcal{F}(\alpha, M, \varepsilon) \subset \mathcal{L}(\alpha, M)$ and $f(x)>\varepsilon$ for every $x$ in the sample space, there exists a randomization $K_{x}$ such that

$$
\sup _{f \in \mathcal{F}(\alpha, M, \varepsilon)}\left\|\overline{\mathbb{P}}_{f} K_{x}-\mathbb{P}_{f}\right\| \leqslant 3 M \varepsilon^{-1 / 2} n^{1 / 2}\left(m^{-\alpha}+m^{-3 / 2}\right) .
$$

The implication is that the experiments are asymptotically equivalent as long as the number
of bins is greater than $n^{-1 /(2 \alpha)}$ for $\alpha<3 / 2$. The transformation in the other direction just bins the continuous observations to produce exactly $\mathbb{\mathbb { P }}_{f}$ from $\mathbb{P}_{f}$. Theorem 2 is proved in Section 2.

Remark. The constants $6 / \sqrt{2 \pi}$ and 3 in these inequalities are not meant to be sharp, but they indicate that the statements are true for a reasonable size constant.

## 2. Binning in density estimation

Let $x_{j}^{*}=(2 j-1) / 2 m$ denote the midpoint of the $j$ th interval. For the $n$ independent observations $X_{1}, \ldots, X_{n}$ with density $f$, let $X_{i}^{*}=x_{j}^{*}$ when $X_{i}$ is in the subinterval $[(j-1) / m, j / m]$. The $X_{1}^{*}, X_{2}^{*}, \ldots, X_{n}^{*}$ are independent observations from a discrete probability distribution $P_{\theta}$ with probabilities

$$
\begin{equation*}
\theta_{j}=\int_{(j-1) / m}^{j / m} f(x) \mathrm{d} x, \quad j=1, \ldots, m \tag{6}
\end{equation*}
$$

on the points $x_{j}^{*}$.
The rounding off of observations onto a regular grid has computational advantages; see Silverman (1982) and Fan and Marron (1994). Hall and Wand (1996) demonstrated that kernel density estimators based on discretized versions of the data perform nearly as well as on the original data. They also discussed 'common linear binning' which is related to our method, but differs in that their triangular kernel is part of the binning procedure while we use the kernel to smooth data that have already been binned.

A sufficient statistic for estimating the $\theta_{j}$ is the number of observations at each $x_{j}^{*}$. Let $\bar{X}_{1}, \ldots, \bar{X}_{m}$ be the counts at each $x_{j}^{*}$. The conditional distribution of the $X_{i}^{*}$ given the $\bar{X}_{i}$ is just the probability of each random ordering of $n$ objects of $m$ types. This conditional distribution does not depend on the $\theta_{j}$, and the necessary randomization takes, for each $j$, $\bar{X}_{j}$ copies of $x_{j}^{*}$ and then randomly assigns an index $i$ to each of the $n$ observations. Therefore, the first part of the randomization $K_{x}$ chooses a random permutation to produce $n$ independent observations from the discrete $P_{\theta}$.

Theorem 2 can then be established by constructing a randomization from the $X_{i}^{*}$ to the $X_{i}$. Working on each coordinate separately, the randomization chooses a new $X_{i}$ to correspond to each $X_{i}^{*}$. By construction, if $X_{i}^{*}=x_{j}^{*}$ then $X_{i}$ must have been in $[(j-1) / m, j / m]$, but we do not know anything else about $X_{i}$. Thus it seems appropriate to have $K_{x}$ uniformly distributed over this subinterval. This results in a continuous marginal distribution that has a piecewise constant density (equal to $m \theta_{j}$ on the $j$ th subinterval). A better approximation to the original distribution of $f$ can be achieved using a $K_{x}$ that is more spread, resulting a smoother density (see Section 4).

Let $K_{j}$ be the conditional distribution of the new $X_{i}$ given that $X_{i}^{*}=x_{j}^{*}$. The marginal distribution $P_{\theta} K_{x}$ is a mixture of $m$ distributions $K_{1}, \ldots, K_{m}$ with weights $\theta_{j}$ such that it has density

$$
\frac{\mathrm{d}\left(P_{\theta} K_{x}\right)}{\mathrm{d} \lambda}(x)=\sum_{j=1}^{m} \theta_{j} \frac{\mathrm{~d} K_{j}}{\mathrm{~d} \lambda}(x) \equiv \hat{f}(x)
$$

where $\lambda$ is uniform on $[0,1]$. Repeating this randomization on the $n$ independent observations defines $K_{x}$ such that $\overline{\mathbb{P}}_{f} K_{x}=P_{\hat{f}}^{n}$. Therefore, the result of the randomization is $n$ independent observations from a distribution with density $\hat{f}$.

The final step is to bound the total variation distance between $n$ independent observations from $f$ and $\hat{f}$, respectively. For product experiments it is worthwhile to use a Hellinger distance bound on the total variation distance

$$
\|\mathbb{P}-\mathbb{Q}\| \leqslant \sqrt{2} H(\mathbb{P}, \mathbb{Q})=\sqrt{2}\left(\frac{1}{2} \int(\sqrt{\mathrm{dP}}-\sqrt{\mathrm{dQ}})^{2}\right)^{1 / 2}
$$

because the squared Hellinger distance between product measures is bounded by the sum of the squared marginal distances,

$$
H^{2}\left(P_{\hat{f}}^{n}, P_{f}^{n}\right) \leqslant n H^{2}\left(P_{\hat{f}}, P_{f}\right)
$$

(see, for example, Strasser 1985, pp. 11-12).
If the densities are bounded away from zero, $f \geqslant \epsilon>0$, then the Hellinger distance can be bounded by the $L^{2}$ distance between the densities,

$$
\frac{1}{2} \int\left(f^{1 / 2}-\hat{f}^{1 / 2}\right)^{2}=\int \frac{(f-\hat{f})^{2}}{2\left(f^{1 / 2}+\hat{f}^{1 / 2}\right)^{2}} \leqslant \frac{1}{2 \epsilon}\|f-\hat{f}\|_{2}^{2}
$$

In Section 4.1 it is shown that for $f \in \mathcal{L}(\alpha, M)$,

$$
\begin{equation*}
\|\hat{f}-f\|_{2}^{2} \leqslant M^{2} m^{-3}+9 M^{2} m^{-2 \alpha} \tag{7}
\end{equation*}
$$

Therefore,

$$
\sup _{f \in \mathcal{F}(\alpha, M, \epsilon)}\left\|P_{\hat{f}}^{n}-\mathbb{P}_{f}^{n}\right\| \leqslant 3 M \epsilon^{-1 / 2} n^{1 / 2}\left(m^{-\alpha}+m^{-3 / 2}\right) .
$$

and Theorem 2 is established.

## 3. Nonparametric regression

The nonparametric regression experiment in (1) has $n$ observations that can be thought of as approximating the increments of the Brownian motion process in (2), $Y_{i} \approx n[Y(i / n)-Y([i-1] / n)]$. Brown and Low (1996) showed that for $g \in \mathcal{L}(\alpha, M)$, $\alpha>1 / 2$, these experiments are equivalent. Using the $Y_{i}$ as approximations to the increments, a continuous Gaussian process can be constructed by interpolating independent Brownian bridges $B_{i}(t)$ on $[(i-1) / n, i / n]$ between the points $Y(i / n)$,

$$
Y_{\bar{g}}(t)=\int_{0}^{t} \sum_{i=1}^{n} Y_{i} \mathbf{1}_{((i-1) / n, i / n]} \mathrm{d} t+\sigma \frac{1}{n} \sum_{i=1}^{n} B_{i}(t) .
$$

This process has mean $\bar{g}(t)=g\left(x_{i}\right)$ for $(i-1) / m \leqslant t<i / m$, and for $0<s<t<1$, $\operatorname{cov}\left(Y_{\bar{g}}(t), Y_{\bar{g}}(s)\right)=\sigma^{2} s / n$. This result can be extended to the unit square if differentiable drift functions can be properly exploited, much as in the density estimation case.

### 3.1. The construction

The constructed process is a function of the $Y_{i}$ and some independent continuous Gaussian processes on $[0,1]^{2}$. These centred processes are determined by their covariance functions (Dudley 2002, Theorem 12.1.3). The Brownian sheet $W(\mathbf{t})$ is a centred Gaussian process with covariance function $C(\mathbf{t}, \mathbf{s})=\left(s_{1} \wedge t_{1}\right)\left(s_{2} \wedge t_{2}\right)$. We also need $K_{i}$-Brownian sheets $W_{K i}$ where the $K_{i}$ are probability measures on $[0,1]^{2}$. Let $\kappa_{i}$ be the cumulative distribution function (cdf) for the measure $K_{i}$. The covariance function for $W_{K i}$ is

$$
\operatorname{cov}\left(W_{K i}(\mathbf{s}), W_{K i}(\mathbf{t})\right)=\kappa_{i}(\mathbf{s} \wedge \mathbf{t}), \quad \mathbf{s}, \mathbf{t} \in[0,1]^{2}
$$

where $\mathbf{s} \wedge \mathbf{t}=\left(\left(s_{1} \wedge t_{1}\right),\left(s_{2} \wedge t_{2}\right)\right)$. Furthermore, a $K_{i}$-Brownian bridge $B_{K i}$ can be constructed via $W_{K i}$ by

$$
B_{K i}(\mathbf{t})=W_{K i}(\mathbf{t})-\kappa_{i}(\mathbf{t}) W_{K i}(1,1) .
$$

The $K_{i}$-Brownian bridge has mean zero and covariance

$$
\begin{equation*}
\operatorname{cov}\left(B_{K i}(\mathbf{s}), B_{K i}(\mathbf{t})\right)=\kappa_{i}(\mathbf{s} \wedge \mathbf{t})-\kappa_{i}(\mathbf{s}) \kappa_{i}(\mathbf{t}) \tag{8}
\end{equation*}
$$

Hence, $\operatorname{var}\left(B_{K i}(\mathbf{t})\right)=\kappa_{i}(\mathbf{t})\left(1-\kappa_{i}(\mathbf{t})\right)$.
The construction of the Gaussian process $Y^{*}(\mathbf{t})$ from the $Y_{i}$ first generates $n$ independent processes $B_{K 1}, \ldots, B_{K n}$, and then

$$
Y^{*}(\mathbf{t})=\frac{1}{n} \sum_{i=1}^{n}\left[Y_{i} \kappa_{i}(\mathbf{t})+\sigma B_{K i}(\mathbf{t})\right] .
$$

The mean of this process is $n^{-1} \sum g\left(\mathbf{x}_{i}^{*}\right) \kappa_{i}(\mathbf{t})$, and the covariance is

$$
\begin{aligned}
\operatorname{cov}\left(Y^{*}(\mathbf{t}), Y^{*}(\mathbf{s})\right) & =\frac{1}{n^{2}} \sum_{i=1}^{n} \sigma^{2} \kappa_{i}(\mathbf{t}) \kappa_{i}(\mathbf{s})+\sigma^{2}\left(\kappa_{i}(\mathbf{s} \wedge \mathbf{t})-\kappa_{i}(\mathbf{s}) \kappa_{i}(\mathbf{t})\right) \\
& =\frac{\sigma^{2}}{n}\left(\frac{1}{n} \sum_{i=1}^{n} \kappa_{i}(\mathbf{s} \wedge \mathbf{t})\right)
\end{aligned}
$$

Assuming an additional condition on the $K_{i}$ measures,

$$
\begin{equation*}
\frac{1}{n} \sum_{i=1}^{n} \kappa_{i}(\mathbf{t})=t_{1} t_{2}, \tag{9}
\end{equation*}
$$

fixes the covariance of $Y_{g}^{*}$ to be the same as $Y_{g}$. The condition in (9) seems reasonable as it implies that if all the $Y_{i}$ are the same then the resulting drift function is constant over the whole square. The kernels described in Section 4.2 meet this criterion, as do uniform distributions over disjoint subsets.

Therefore, $Y^{*}$ actually has the same distribution as

$$
Y_{\hat{g}}(\mathbf{t})=\int_{0}^{t_{1}} \int_{0}^{t_{2}} \hat{g}(\mathbf{x}) \mathrm{d} \mathbf{x}+\frac{\sigma}{\sqrt{n}} W(\mathbf{t})
$$

where

$$
\hat{g}(\mathbf{x})=\frac{1}{n} \sum_{i=1}^{n} g\left(\mathbf{x}_{i}^{*}\right) \frac{\mathrm{d} K_{i}}{\mathrm{~d} \lambda}(\mathbf{x}) .
$$

This should be close in distribution to the $Y_{g}$ if $g$ and $\hat{g}$ are close.

### 3.2. A total variation distance bound

Shifting two distributions by the same amount will not affect the total variation distance between them. Thus, without loss of generality, we will find the total variation distance between a process $Y_{0}(\mathbf{t})$ with mean 0 and variance $\sigma^{2} / n$, and a process $Y_{\delta}(\mathbf{t})$ with drift $\delta(\mathbf{t})$ and the same variance. Let $\mathbb{Q}_{0}$ and $\mathbb{Q}_{\delta}$ be the distributions of $Y_{0}$ and $Y_{\delta}$, respectively. The set of sample paths that achieves the supremum in (3) is the set where

$$
\frac{\mathrm{d} \mathbb{Q}_{\delta}}{\mathrm{d} \mathbb{Q}_{0}}=\exp \left[\frac{n}{\sigma^{2}}\left(\int \delta(\mathbf{t}) \mathrm{d} Y(\mathbf{t})-\frac{1}{2}\|\delta\|_{2}^{2}\right)\right]>1
$$

Let $A$ be this set of continuous sample paths such that $\int \delta(\mathbf{t}) \mathrm{d} Y(\mathbf{t})>\frac{1}{2}\|\delta\|_{2}^{2}$. Under $\mathbb{Q}_{0}$, the integral $\int \delta(\mathbf{t}) \mathrm{d} Y(\mathbf{t})$ has a normal distribution with mean 0 and variance $\sigma^{2}\|\delta\|_{2}^{2} / n$. Under $\mathbb{Q}_{\delta}$, the integral has mean $\|\delta\|_{2}^{2}$ and the same variance. Therefore the total variation distance is

$$
\begin{equation*}
\left\|\mathbb{Q}_{0}-\mathbb{Q}_{\delta}\right\|=\left|\mathbb{Q}_{0} A-\mathbb{Q}_{\delta} A\right|=1-2 \Phi(-\Delta / 2) \tag{10}
\end{equation*}
$$

where $\Delta=\sigma^{-1} n^{1 / 2}\|\delta\|_{2}$. The expression in (10) for the total variation distance is concave for positive $\Delta$, so a simple expansion gives

$$
\left\|\mathbb{Q}_{0}-\mathbb{Q}_{\delta}\right\| \leqslant \frac{1}{\sqrt{2 \pi}} \Delta .
$$

In the case of $Y_{g}$ and $Y_{\hat{g}}$, the bound again depends essentially on the $L_{2}$ distance between the means,

$$
\left\|\mathbb{Q}_{g}-\mathbb{Q}_{\hat{g}}\right\| \leqslant \frac{1}{\sqrt{2 \pi}} \frac{\sqrt{n}}{\sigma}\|g-\hat{g}\|_{2} .
$$

A bound on the $L^{2}$ distance is needed, in Section 4.2 we will show that

$$
\begin{equation*}
\sup _{g \in \mathcal{L}(\alpha, M)}\|g-\hat{g}\|_{2} \leqslant 6 M\left(n^{-\alpha / 2}+n^{-3 / 4}\right) . \tag{11}
\end{equation*}
$$

Therefore,

$$
\sup _{g \in \mathcal{L}(\alpha, M)}\left\|\mathbb{Q}_{g} K_{y}-\mathbb{Q}_{g}\right\| \leqslant \frac{6 M}{\sqrt{2 \pi \sigma^{2}}}\left(n^{1 / 2-\alpha / 2}+n^{-1 / 4}\right),
$$

which proves Theorem 1.

## 4. Bounding the $L_{2}$ distances

The $K_{j}$ are probability measures with densities equal to linear interpolating functions:

where $x_{j-1}^{*}=x_{j}^{*}-1 / m$ and $x_{j+1}^{*}=x_{j}^{*}+1 / m$.
For $j=1$, we will need a different conditional density:

to avoid getting observations outside [0, 1]. The analogous measure will be used for $j=m$ to avoid results greater than 1. Equivalently, we could use the triangular measures everywhere and then reflect any observations outside the interval back in.

The average

$$
\frac{1}{m} \sum_{j=1}^{m} a_{j} \frac{\mathrm{~d} K_{j}}{\mathrm{~d} \lambda}(x)
$$

is a piecewise linear function that is equal to $a_{j}$ at the midpoints $x_{j}^{*}$.

### 4.1. One dimension

For the set of functions $f \in \mathcal{L}(\alpha, M)$ with continuous first derivatives ( $\alpha>1$ ), the error in a linear Taylor expansion is

$$
\begin{equation*}
\left|f(t+\delta)-f(t)-\delta f^{\prime}(t)\right| \mathrm{h}=\left|f(t)+\delta f^{\prime}\left(t^{*}\right)-f(t)-\delta f^{\prime}(t)\right| \leqslant M \delta^{\alpha} \tag{12}
\end{equation*}
$$

from (4). The existence of $t^{*}$ such that $\left|t^{*}-t\right| \leqslant \delta$ and $f(t+\delta)=f(t)+\delta f^{\prime}\left(t^{*}\right)$ follows from the mean value theorem.

Any $x$ between $(2 m)^{-1}$ and $1-(2 m)^{-1}$ lies between two of the grid points $x_{j}^{*}$ and $x_{j+1}^{*}$. The difference between $f$ and $\hat{f}$ at such an $x$ is

$$
\begin{equation*}
|f(x)-\hat{f}(x)| \leqslant\left|f\left(x_{j}^{*}\right)-\hat{f}\left(x_{j}^{*}\right)\right|+\left|x-x_{j}^{*}\right|\left|f^{\prime}\left(x_{j}^{*}\right)-\hat{f}^{\prime}\left(x_{j}^{*}\right)\right|+\left|E_{1}\right|+\left|E_{2}\right|, \tag{13}
\end{equation*}
$$

where $E_{1}$ and $E_{2}$ are the errors in expansions around $x_{j}^{*}$ of $f$ and $\hat{f}$, respectively. By (12), $\left|E_{1}\right| \leqslant M m^{-\alpha}$.

Let us bound the first term on the right in (13). The average over any interval $[(i-1) / m, i / m]$ can be approximated using an expansion around the midpoint $x_{j}^{*}$,

$$
\begin{aligned}
\hat{f}\left(x_{j}^{*}\right) & =m \int_{(i-1) / m}^{i / m} f(x) \mathrm{d} x \\
& =m \int_{(i-1) / m}^{i / m} f\left(x_{j}^{*}\right)+\left(x-x_{j}^{*}\right) f^{\prime}\left(x_{j}^{*}\right)+E_{3} \mathrm{~d} x \\
& =f\left(x_{j}^{*}\right)+E_{4}
\end{aligned}
$$

because $\int_{(i-1) / m}^{i / m}\left(x-x_{j}^{*}\right)=0$. The bound in (12) implies that $\left|E_{3}\right| \leqslant M m^{-\alpha}$ and thus $\left|E_{4}\right| \leqslant M m^{-\alpha}$.

The second term on the right in (13) is problematic because the derivative $\hat{f}^{\prime}\left(x_{j}^{*}\right)$ does not exist. However, if we are only interested in $x_{j}^{*}<x<x_{j+1}^{*}$ then defining $\hat{f}^{\prime}\left(x_{j}^{*}\right)=m\left(\hat{f}\left(x_{j}^{*}\right)-\hat{f}\left(x_{j+1}^{*}\right)\right)$ makes $\left|E_{2}\right| \equiv 0$. Making two appeals to the mean value theorem,

$$
\begin{aligned}
\hat{f}^{\prime}\left(x_{j}^{*}\right) & =m\left(\hat{f}\left(x_{j}^{*}\right)-\hat{f}\left(x_{j+1}^{*}\right)\right) \\
& =m^{2} \int_{(i-1) / m}^{i / m} f(x)-f\left(x+m^{-1}\right) \mathrm{d} x \\
& =m^{2} \int_{(i-1) / m}^{i / m} m^{-1} f^{\prime}\left(\xi_{x}\right) \mathrm{d} x \\
& =f^{\prime}\left(\xi_{j}\right),
\end{aligned}
$$

where $\xi_{x}$ and $\xi_{j}$ are arbitrary points in the interval $\left[x_{j-1}^{*}, x_{j+1}^{*}\right]$ that make the statements true. Therefore, the second term is

$$
\left|x-x_{j}^{*}\right|\left|f^{\prime}\left(x_{j}^{*}\right)-\hat{f}^{\prime}\left(x_{j}^{*}\right)\right| \leqslant \frac{1}{m}\left|f^{\prime}\left(x_{j}^{*}\right)-f^{\prime}\left(\xi_{j}\right)\right| \leqslant M m^{-\alpha},
$$

and each of the terms in (13) is bounded:

$$
\begin{equation*}
|f(x)-\hat{f}(x)| \leqslant\left|E_{4}\right|+M m^{-\alpha}+\left|E_{1}\right|+0 \leqslant 3 M m^{-\alpha} \tag{14}
\end{equation*}
$$

This argument does not work at the edges. If $x<(2 m)^{-1}$ then $\hat{f}(x)=\hat{f}\left([2 m]^{-1}\right)$ and

$$
\begin{align*}
|f(x)-\hat{f}(x)| & \leqslant\left|\hat{f}\left([2 m]^{-1}\right)-f\left([2 m]^{-1}\right)\right|+\left|x-x_{j}^{*}\right|\left|f^{\prime}\left(x_{j}^{*}\right)\right|+\left|E_{1}\right| \\
& \leqslant 2 M m^{-\alpha}+\frac{1}{2} M m^{-1} \tag{15}
\end{align*}
$$

because $\left|f^{\prime}\left(x_{j}^{*}\right)\right|<M$ by assumption. The same argument works for $x>1-(2 m)^{-1}$.
Squaring the pointwise bounds in (14) and (15) and then integrating,

$$
\begin{align*}
\|\hat{f}-f\|_{2}^{2} & =\int_{0}^{1 /(2 m)}(\hat{f}-f)^{2}+\int_{1 /(2 m)}^{1-1 /(2 m)}(\hat{f}-f)^{2}+\int_{1-1 /(2 m)}^{1}(\hat{f}-f)^{2} \\
& \leqslant m^{-1} M^{2} m^{-2}+9 M^{2} m^{-2 \alpha} \tag{16}
\end{align*}
$$

Therefore, the bound needed in (7) is established.

### 4.2. Two dimensions

For regression mean functions $g(\mathbf{x})$ on $[0,1]^{2}$, we will use kernels that are products of the one-dimensional kernels above. The cdf of the kernel $K_{i, j}$ is $\kappa_{i, j}(\mathbf{x})=\kappa_{i}\left(x_{1}\right) \kappa_{j}\left(x_{2}\right)$, where $\kappa_{i}$ and $\kappa_{j}$ are the cdfs of $K_{i}$ and $K_{j}$ respectively.

In $\mathbb{R}^{2}$, the differentiability of $g$ gives

$$
\begin{equation*}
g(\mathbf{x})=g\left(\mathbf{x}_{i, j}^{*}\right)+\left(\mathbf{x}-\mathbf{x}_{i, j}^{*}\right)^{\mathrm{T}} g^{\prime}\left(\mathbf{x}_{i, j}^{*}\right)+E_{1}, \tag{17}
\end{equation*}
$$

and the Lipschitz condition implies that the error is bounded by

$$
\begin{equation*}
\left|E_{1}\right| \leqslant\left|\mathbf{x}-\mathbf{x}_{i, j}^{*}\right|\left|g^{\prime}\left(\mathbf{x}_{i, j}^{*}\right)-g^{\prime}(\xi)\right| \leqslant M\left|\mathbf{x}-\mathbf{x}_{i, j}^{*}\right|^{\alpha} \tag{18}
\end{equation*}
$$

as in the one-dimensional case.
The mean value of the constructed process $\hat{g}$ is no longer linear, but the quadratic part is small. Consider points $\mathbf{x}$ in the interior of the square formed by the four midpoints $\mathbf{x}_{i, j}^{*}, \mathbf{x}_{i+1, j}^{*}, \mathbf{x}_{i, j+1}^{*}$, and $\mathbf{x}_{i+1, j+1}^{*}$. We have

$$
\begin{equation*}
\hat{g}(\mathbf{x})=g\left(\mathbf{x}_{i, j}^{*}\right)+\sqrt{n}\left(\mathbf{x}-\mathbf{x}_{i, j}^{*}\right)^{\mathrm{T}}\binom{g\left(\mathbf{x}_{i+1, j}^{*}\right)-g\left(\mathbf{x}_{i, j}^{*}\right)}{g\left(\mathbf{x}_{i, j+1}^{*}\right)-g\left(\mathbf{x}_{i, j}^{*}\right)}+E_{2} . \tag{19}
\end{equation*}
$$

The differences can be written as

$$
\begin{aligned}
& g\left(\mathbf{x}_{i+1, j}^{*}\right)-g\left(\mathbf{x}_{i, j}^{*}\right)=\left(\begin{array}{ll}
\frac{1}{\sqrt{n}} & 0
\end{array}\right) g^{\prime}\left(\mathbf{x}_{i, j}^{*}\right)+E_{3}, \\
& g\left(\mathbf{x}_{i, j+1}^{*}\right)-g\left(\mathbf{x}_{i, j}^{*}\right)=\left(\begin{array}{ll}
0 & \frac{1}{\sqrt{n}}
\end{array}\right) g^{\prime}\left(\mathbf{x}_{i, j}^{*}\right)+E_{4} .
\end{aligned}
$$

The bound in (18) means the errors $\left|E_{3}\right|$ and $\left|E_{4}\right|$ are less than $M n^{-\alpha / 2}$.
The error term $E_{2}$ in (19) is the quadratic component of the mean function. Let $\left(\zeta_{1}, \zeta_{2}\right)=\left(\mathbf{x}-\mathbf{x}_{i, j}^{*}\right)^{\mathrm{T}}$; then

$$
E_{2}=n \zeta_{1} \zeta_{2}\left[g\left(\mathbf{x}_{i, j}^{*}\right)-g\left(\mathbf{x}_{i+1, j}^{*}\right)-g\left(\mathbf{x}_{i, j+1}^{*}\right)+g\left(\mathbf{x}_{i, j+1}^{*}\right)\right] .
$$

The size of $\left|E_{2}\right|$ is bounded using $n\left|\zeta_{1} \zeta_{2}\right| \leqslant 1$ and

$$
\left|g\left(\mathbf{x}_{i, j}^{*}\right)-g\left(\mathbf{x}_{i+1, j}^{*}\right)-g\left(\mathbf{x}_{i, j+1}^{*}\right)+g\left(\mathbf{x}_{i+1, j+1}^{*}\right)\right| \leqslant n^{-1 / 2}\left|g^{\prime}\left(\mathbf{x}_{i, j+1}^{*}\right)-g^{\prime}\left(\mathbf{x}_{i, j}^{*}\right)\right|+\left|E_{3}\right|+\left|E_{5}\right|,
$$

where $E_{5}$ is the error in approximating the difference $g\left(\mathbf{x}_{i+1, j+1}^{*}\right)-g\left(\mathbf{x}_{i, j+1}^{*}\right)$ by the derivative at $\mathbf{x}_{i, j+1}^{*}$. Thus, $\left|E_{5}\right| \leqslant M n^{-\alpha / 2}$ and $\left|E_{2}\right| \leqslant 3 M n^{-\alpha / 2}$.

Putting together (17) and (19),

$$
\begin{aligned}
|g(\mathbf{x})-\hat{g}(\mathbf{x})| \leqslant & \mid g\left(\mathbf{x}_{i, j}^{*}\right)+\left(\mathbf{x}-\mathbf{x}_{i, j}^{*}\right)^{\mathrm{T}} g^{\prime}\left(\mathbf{x}_{i, j}^{*}\right)+E_{1} \\
& -\left[g\left(\mathbf{x}_{i, j}^{*}\right)+\left(\mathbf{x}-\mathbf{x}_{i, j}^{*}\right)^{\mathrm{T}} g^{\prime}\left(\mathbf{x}_{i, j}^{*}\right)+n^{1 / 2} \zeta_{1} E_{3}+n^{1 / 2} \zeta_{2} E_{4}+E_{2}\right] \mid \leqslant 6 M n^{-\alpha / 2}
\end{aligned}
$$

for $\mathbf{x}$ not within $\frac{1}{2} n^{-1 / 2}$ of the edge of the square. The contribution from the centre of the square to $\|g-\hat{g}\|_{2}^{2}$ is therefore less than $36 M^{2} n^{-\alpha}$.

To see what happens near the edges of the square, consider a point $\mathbf{x}$ where the first coordinate is less than $\frac{1}{2} n^{-1 / 2}$; then the $\hat{g}$ function is exactly

$$
\hat{g}(\mathbf{x})=g\left(\mathbf{x}_{i, j}^{*}\right)+\sqrt{n}\left(\mathbf{x}-\mathbf{x}_{i, j}^{*}\right)^{\mathrm{T}}\binom{g\left(\mathbf{x}_{i+1, j}^{*}\right)-g\left(\mathbf{x}_{i, j}^{*}\right)}{0} .
$$

Thus,

$$
|g(\mathbf{x})-\hat{g}(\mathbf{x})| \leqslant\left|E_{1}\right|+\left|E_{3}\right|+\frac{M}{2} n^{-1 / 2} \leqslant \frac{1}{2} M n^{-1 / 2}+\frac{3}{2} M n^{-\alpha / 2}
$$

where $\frac{1}{2} \mathrm{Mn}^{-1 / 2}$ bounds the contribution from the second coordinate of $g^{\prime}(\mathbf{x})$ because $\left|g^{\prime}(\mathbf{x})\right| \leqslant M$. An analogous bound can be put on the errors along the other sides.

Finally, if $x_{1}<\frac{1}{2} n^{-1 / 2}$ and $x_{2}<\frac{1}{2} n^{-1 / 2}$, then $\hat{g}=g\left(\mathbf{x}_{i, j}^{*}\right)$ and thus $\left|g(\mathbf{x})-g\left(\mathbf{x}_{i, j}^{*}\right)\right|$ $\leqslant M(2 n)^{-1 / 2}$. The same bound applies in the other three corners.

The total contribution to $\|g-\hat{g}\|_{2}^{2}$ from all the area near the edges is less than $M^{2} n^{-3 / 2} / 2$ because the total area is less than $2 n^{-1 / 2}$.

Therefore,

$$
\|g-\hat{g}\|_{2}^{2} \leqslant \frac{1}{2} M^{2} n^{-3 / 2}+36 M^{2} n^{-\alpha}
$$

which establishes (11).

### 4.3. The other direction

Asymptotic equivalence also requires a transformation in the other direction: a way to generate the $n$ regression observations from the continuous Gaussian process. The transformation uses the increments of the process over each square,

$$
Y_{i, j}^{*}=n\left[Y\left(\frac{i}{\sqrt{n}}, \frac{j}{\sqrt{n}}\right)-Y\left(\frac{i-1}{\sqrt{n}}, \frac{j}{\sqrt{n}}\right)-Y\left(\frac{i}{\sqrt{n}}, \frac{j-1}{\sqrt{n}}\right)+Y\left(\frac{i-1}{\sqrt{n}}, \frac{j-1}{\sqrt{n}}\right)\right] .
$$

Let $I_{i, j}=\left\{\left(x_{1}, x_{2}\right):(i-1) / \sqrt{n}<x_{1} \leqslant i / \sqrt{n},(j-1) / \sqrt{n}<x_{2} \leqslant j / \sqrt{n}\right\}$. Then the increments $Y_{i, j}^{*}$ are independent with variance $\sigma$ and have mean $n \iint_{i, j} g \mathrm{~d} \mathbf{x}$.

The only error then is the difference between $g\left(\mathbf{x}_{i, j}^{*}\right)$ and this average:

$$
\begin{aligned}
g\left(\mathbf{x}_{i, j}^{*}\right)-n \iint_{I_{i, j}} g(\mathbf{x}) \mathrm{d} \mathbf{x} & =n \iint_{I_{i, j}} g\left(\mathbf{x}_{i, j}^{*}\right)-g(\mathbf{x}) \mathrm{d} \mathbf{x} \\
& =n \iint_{I_{i, j}}\left(\mathbf{x}-\mathbf{x}_{i, j}^{*}\right)^{\mathrm{T}} g^{\prime}\left(\mathbf{x}_{i, j}^{*}\right)+E_{1} \mathrm{~d} \mathbf{x} \\
& =n \iint_{I_{i, j}} E_{1} \mathrm{~d} \mathbf{x} .
\end{aligned}
$$

Therefore, by (18), $\left(g\left(\mathbf{x}_{i, j}^{*}\right)-n \iint_{I_{i, j}} g(\mathbf{x}) \mathrm{d} \mathbf{x}\right)^{2} \leqslant M n^{-\alpha}$. The total variation distance between the joint distributions of all $n$ observations is therefore less than $(\sqrt{2 \pi} \sigma)^{-1} M n^{-\alpha / 2+1 / 2}$.

This error is less than that made in the other direction and therefore whenever the bound in Theorem 1 goes to 0 , the experiments are asymptotically equivalent.

### 4.4. Higher dimensions

The same technique could provide a solution in the three-dimensional case if there is an added restriction to the class of parameter functions in order to lessen the edge effects.

The kernel would still be a product of the one-dimensional kernels in each of the dimensions $K(\mathbf{x})=K_{i}\left(x_{1}\right) K_{j}\left(x_{2}\right) K_{k}\left(x_{3}\right)$. The errors in the inner part of the cube are of order $O\left(n^{-\alpha / 3}\right)$.

The points near the edge of the cube, however, would make an error of order $n^{-1 / 3}$ in at least one of the three coordinates. This contributes a term of order $O\left(n^{-1}\right)$ to the squared difference between $g$ and $\hat{g}$. The result is that the total variation distance between the Gaussian processes will not converge to 0 . It is possible to impose conditions on the drift functions to minimize these edge effects. For example, imposing a periodic boundary condition and adjusting the edge kernels to be periodic makes the error $O\left(n^{-\alpha / 3}\right)$ everywhere. This is sufficient for asymptotic equivalence if $\alpha>3 / 2$.

Higher dimensions would require $\alpha>2$ (see Brown and Zhang 1998), and our methods cannot take advantage of this added smoothness. Higher-order interpolating kernels are negative in places, and the requirement that the kernels have a positive density function is critical in the construction.

Using a Fourier series expansion, Rohde (2004) was able to take advantage of arbitrary amounts of smoothness in the case of periodic functions in the nonparametric regression problem, and it is likely those techniques could be extended to higher-dimensional problems.

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